

On asymptotic behavior of Heine-Stieltjes and Van Vleck polynomials

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This paper is dedicated to the 60th Birthday of Guillermo (Bill) López Lagomasino

ABSTRACT. We investigate the strong asymptotics of Heine-Stieltjes polynomials – polynomial solutions of a second order differential equations with complex polynomial coefficients. The solution is given in terms of critical measures (saddle points of the weighted logarithmic energy on the plane), that are tightly related to quadratic differentials with closed trajectories on the plane. The paper is a continuation of the research initiated in [8]. However, the starting point here is the WKB method, which allows to obtain the strong asymptotics.

1. Generalized Lamé equation

In 1878 Heine studied the following problem: given two polynomials,

$$(1.1) \quad A(z) = \prod_{i=0}^p (z - a_i), \quad B(z) = \alpha z^p + \text{lower degree terms}, \quad \alpha \in \mathbb{C},$$

where the zeros of A are complex and pairwise distinct, describe the polynomial solutions of the *generalized Lamé differential equation* (in algebraic form),

$$(1.2) \quad A(z) y''(z) + B(z) y'(z) - n(n + \alpha - 1) V_n(z) y(z) = 0,$$

where V_n is a polynomial (in general, depending on n) of degree $\leq p - 1$; if $\deg V = p - 1$, then V is monic. For $p = 1$ we can easily recognize in (1.2) the hypergeometric differential equation, while for $p = 2$ it is known as the Heun's equation (see [15]).

Heine [4] proved that for every $n \in \mathbb{N}$ there exist at most

$$(1.3) \quad \sigma(n) = \binom{n + p - 1}{n}$$

different polynomials V_n such that (1.2) admits a polynomial solution $y = Q_n \in \mathbb{P}_n$; hereafter \mathbb{P}_n stands for the set of all algebraic polynomials of degree $\leq n$. These

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particular V_n 's are called *Van Vleck polynomials*, and the corresponding polynomial solutions $y = Q_n$ are known as *Heine-Stieltjes* (or simply Stieltjes) polynomials. Furthermore, if the polynomials A and B are algebraically independent then for any $n \in \mathbb{N}$ there exist *exactly* $\sigma(n)$ Van Vleck polynomials V_n , their degree is *exactly* $p-1$, and for each V_n equation (1.2) has a unique (up to a constant factor) solution y of degree n . According to recent results in [16], this statement still holds for arbitrary A and B for all sufficiently large n .

Stieltjes discovered an electrostatic interpretation of zeros of the polynomials discussed in [4], which attracted common attention to the problem. He studied the problem (1.2) in a particular setting, assuming that all zeros of A are real and that all residues ρ_k in

$$(1.4) \quad \frac{B(x)}{A(x)} = \sum_{k=0}^p \frac{\rho_k}{x - a_k}$$

are strictly positive. He proved in [18] (see also [20, Theorem 6.8]) that in this case for each $n \in \mathbb{N}$ there are *exactly* $\sigma(n)$ different Van Vleck polynomials of degree $p-1$ and the same number of corresponding Heine-Stieltjes polynomials y of degree n , given by all possible ways how the n zeros of y can be distributed in the p open intervals defined by the zeros of A .

After these two contributions a vast number of research papers on this topic has been published, dealing mainly with the real situation considered by Stieltjes. Besides the relevant contribution to the algebraic theory [16], we can mention here the paper [9], where the limit distribution of zeros of Heine-Stieltjes polynomials for the Stieltjes case has been established in terms of the traditional extremal problem for the weighted logarithmic energy on a compact set of the plane.

A new approach to the asymptotics of the polynomial solutions of (1.2), based on a different type of equilibrium problem, has been developed in [8], where zero distribution of Heine-Stieltjes polynomials was investigated. The present paper is in many senses a continuation of [8]. We use essentially the same tools: families of continuous critical measures and their representations in terms of quadratic differentials with closed trajectories. However, the starting point of our considerations here is different: instead of electrostatics we use the WKB method, which allows to obtain the strong asymptotics.

2. Asymptotic formula for Heine-Stieltjes polynomials

In this section we describe briefly the strong asymptotics of polynomial solutions of (1.2). As we will see immediately, it is based on the asymptotics of Van Vleck polynomials, so formally we should have started with the latter topic. We choose another way of presentation, which follows rather the logic of the proof of the main result. We formulate next a theorem on asymptotics of Heine-Stieltjes polynomials, which contains Van Vleck polynomials as unknown (“access”) parameters.

Hereafter we assume that $n \in \mathbb{N}$ is sufficiently large, so that there are exactly $\sigma(n)$ Van Vleck polynomials, each of degree exactly $p-1$, see [16]; denote by V_n one of these Van Vleck polynomials, and let \mathcal{V}_n stand for the set of zeros of V_n , $\mathcal{A} = \{a_0, \dots, a_p\}$ the set of zeros of A , and $\Omega_n \stackrel{\text{def}}{=} \mathbb{C} \setminus (\mathcal{A} \cup \mathcal{V}_n)$.

Fix any $z_0 \in \mathbb{C}$; in any simply-connected domain $D \subset \Omega$ we can select a single-valued branch of $\sqrt{V_n/A}$ and define the *natural parameter*

$$(2.1) \quad \xi = \xi_n(z) = \int_{z_0}^z \sqrt{\frac{V_n(t)}{A(t)}} dt.$$

In what follows, many statements will be made in terms of trajectories of the quadratic differential

$$\varpi_n = -\frac{V_n(z)}{A(z)} (dz)^2,$$

which are basically level lines of the function $\text{Im } \xi_n$.

We define also

$$(2.2) \quad \zeta_n(z) = \exp(\xi_n(z)) = \exp\left(\int_{z_0}^z \sqrt{\frac{V_n(t)}{A(t)}} dt\right),$$

as well as the parameter

$$(2.3) \quad \lambda_n = n + \frac{\alpha - 1}{2}.$$

Theorem 2.1. *With the assumptions and notations above, let $y = Q_n$ be a Heine-Stieltjes polynomial (solution of (1.2)) corresponding to V_n . Then:*

- (i) *there exist a set $\Gamma_n = \cup_{k \geq 1} \gamma_{k,n}$ comprised of a finite number of at most $2p - 1$ disjoint Jordan arcs $\gamma_{k,n}$, such that for an appropriate selection of the branch in (2.1),*

$$(2.4) \quad y(z) = H_n(z) \zeta_n^{\lambda_n}(z) (1 + \varepsilon_n(z)), \quad z \in \mathbb{C} \setminus \Gamma_n,$$

where

$$(2.5) \quad H_n(z) = \left(\frac{A}{V_n}(z)\right)^{1/4} \exp\left(-\int^z \frac{B}{2A}(t) dt\right).$$

Moreover, $\varepsilon_n(z) \rightarrow 0$ with $n \rightarrow \infty$ uniformly for $\text{dist}(z, \Gamma_n) \geq C > 0$.

- (ii) *The endpoints of each $\gamma_{k,n}$ belong to $\mathcal{A} \cup \mathcal{V}_n$. Two arcs $\gamma_{j,n}$ and $\gamma_{k,n}$ for $i \neq j$ are either disjoint or have a common endpoint. The domain $\overline{\mathbb{C}} \setminus \Gamma_n$ is connected, and $\sqrt{V_n/A}$ has a single-valued branch. Consequently, the branch of the square root in (2.4) is fully determined by condition*

$$\lim_{z \rightarrow \infty} z \sqrt{\frac{V_n(z)}{A(z)}} = 1.$$

- (iii) *Each $\gamma_{k,n}$ can be selected as a part of a trajectory of the quadratic differential*

$$(2.6) \quad \varpi_n = -\frac{V_n(z)}{A(z)} (dz)^2,$$

which is close to its critical trajectory completed by two small segments connecting the trajectory with two points from $\mathcal{A} \cup \mathcal{V}_n$.

Observe that formula (2.4) is given in terms of polynomials V_n . In this way, the problem of asymptotics of Heine-Stieltjes polynomials is reduced to that of Van Vleck polynomials. We discuss the location of the zeros of V_n 's in Section 4.

The critical role is played by the set $\Gamma_n = \cup_{k=1}^p \gamma_{k,n}$. More exactly, formula (2.4) should be combined with the local asymptotic formula for Q_n in a neighborhood of $\gamma_{k,n}$ (but away from the endpoints):

$$(2.7) \quad Q_n(z) = H_n(z) \left[\zeta_n^{\lambda_n}(z) (1 + \varepsilon_{n,1}(z)) + \zeta_n^{-\lambda_n}(z) (1 + \varepsilon_{n,2}(z)) \right],$$

for certain small $\varepsilon_{n,j}(z)$. Local formulas of this type can be obtained directly from the WKB estimates as explained in [11]. However, the existence of a global formula with the specified dominant term is not totally trivial. An arc $\gamma_{k,n}$ enters the WKB analysis naturally as the Stokes line: the dominant term in (2.7) changes across $\gamma_{k,n}$ (Stokes' phenomenon). This means that two terms are asymptotically equal on $\gamma_{k,n}$, and equation

$$\zeta_n^{2\lambda_n}(z) = -1$$

defines the zeros of Q_n . In other words, arcs $\gamma_{k,n}$ asymptotically carry zeros of Q_n .

The proof of the theorem is based on the Liouville-Green (also known as the WKB) asymptotic formula. As we have mentioned above, the WKB method, the way it is explained in [11], is ready to use in order to obtain local estimates. Construction of a global formula from the local ones constitutes a challenge, as it usually happens in the WKB analysis. So, the proof of the theorem, which is outlined next, is not entirely on the surface, and relies heavily on the fact that both the coefficients and a solution of (1.2) are polynomials.

3. WKB analysis of the Heine-Stieltjes polynomials

The Liouville-Green (also known as the WKB) method is a standard tool of asymptotic analysis of ordinary differential equations; we refer the reader to the comprehensive account contained in the monographs [1] and [11].

It was probably Nuttall [10] who applied the WKB analysis to the study of the strong asymptotics of the generalized Jacobi polynomials via a differential equation of the form (1.2).

3.1. Liouville's transformation. Recall that at this stage we assume that the Van Vleck polynomial V_n is fixed, and denote by $y = Q_n$ a corresponding Heine-Stieltjes polynomial of degree n . The following procedure is well known, so we summarize it here very briefly.

Rewriting equation (1.2) in terms of

$$(3.1) \quad u(z) = y(z) \exp \left(\int^z \frac{B}{2A}(t) dt \right),$$

we eliminate the first order derivative and get $u''(z) = F_n(z)u(z)$, with

$$(3.2) \quad F_n(z) = (\lambda_n^2 - \rho^2) \frac{V_n}{A}(z) + \left(\frac{B}{2A} \right)^2 (z) + \left(\frac{B}{2A} \right)' (z),$$

where λ_n has been defined in (2.3) and

$$\rho = \frac{\alpha - 1}{2}.$$

Fix any $z_0 \in \mathbb{C}$; the Liouville-Green transformation $w(\xi) = (\xi')^{1/2} u$, with $\xi = \xi_n(z)$ defined in (2.1), yields

$$(3.3) \quad \frac{d^2 w}{d\xi^2} = (\lambda_n^2 + g_n(\xi)) w(\xi),$$

where

$$(3.4) \quad g_n(\xi) = G_n(z(\xi)), \quad G_n = \frac{A}{V_n} \left(\left(\frac{B}{2A} \right)^2 + \left(\frac{B}{2A} \right)' - T_n^2 - T_n' \right) - \rho^2,$$

and

$$T_n(z) = \frac{1}{4} \left(\frac{A'}{A} - \frac{V_n'}{V_n} \right) (z).$$

A detailed discussion of this transformation in the complex plane and of the resulting WKB approximation can be found, e.g. in [11].

3.2. Quadratic differentials on the Riemann sphere. At this point let us recall briefly some notions from the theory of rational quadratic differentials on the Riemann sphere; see [13] and [19] for further details. Given the quadratic differential ϖ_n defined in (2.6), a smooth curve γ along which

$$\varpi_n = -V_n(z)/A(z) (dz)^2 > 0 \quad \Leftrightarrow \quad \operatorname{Im} \xi_n(z) = \operatorname{const}$$

is a *horizontal arc* of ϖ_n ; the distinguished parameter ξ_n was defined in (2.1). More precisely, if γ is given by a parametrization $z(t)$, $t \in (\alpha, \beta)$, then

$$-\frac{V_n}{A}(z(t)) \left(\frac{dz}{dt} \right)^2 > 0, \quad t \in (\alpha, \beta).$$

A maximal horizontal arc is called a *horizontal trajectory* (or simply a *trajectory*) of ϖ_n . Analogously, trajectories of $-\varpi_n$ are called *orthogonal* or *vertical trajectories* of ϖ_n ; along these curves

$$V_n(z)/A(z) (dz)^2 > 0 \quad \Leftrightarrow \quad \operatorname{Re} \xi_n(z) = \operatorname{const}.$$

We can define a conformal invariant metric associated with the quadratic differential ϖ , given by the length element $|d\xi_n| = |\sqrt{V_n/A}|(z)|dz|$; the ϖ_n -length of a curve γ is

$$\|\gamma\|_{\varpi_n} = \frac{1}{\pi} \int_{\gamma} \sqrt{\left| \frac{V_n}{A} \right|} (z) |dz|;$$

(observe that this definition differs by a normalization constant from the definition 5.3 in [19]). Furthermore, if D is a simply connected domain not containing singular points of ϖ_n , we can introduce the ϖ_n -distance by

$$\operatorname{dist}(z_1, z_2; \varpi, D) = \inf \{ \|\gamma\|_{\varpi_n} : z_1, z_2 \in \bar{\gamma}, \gamma \subset D \}.$$

Trajectories and orthogonal trajectories are in fact geodesics (in the ϖ_n -metric) connecting any two of its points (see [13, Thm. 8.4]).

A simply connected domain D not containing points from $\mathcal{A} \cup \mathcal{V}_n$ is called a ϖ_n -*rectangle* if it is delimited by two horizontal and two vertical arcs of ϖ_n ; in other words, if $\xi_n(D)$ is a (euclidean) rectangle $[a, b] \times [c, d]$, and $D \mapsto \xi_n(D)$ is a one-to-one conformal mapping. We call the value $d - c$ the ϖ_n -height, $b - a$ the ϖ_n -length, and $2(b - a + d - c)$ the ϖ_n -perimeter of D . Obviously, these definitions are consistent with the freedom in the selection of the natural parameter ξ_n .

3.3. Local asymptotics. Theorems on local asymptotics of solutions are well-known (see [1] and [11]); using the local estimates that appear in [11, Ch. VI, Theorem 11.1] it follows that if D is a ϖ_n -rectangle such that function g_n defined in (3.4) is holomorphic in $\xi_n(D) = [a, b] \times [c, d]$, then the differential equation (3.3) has in $[a, b] \times [c, d]$ two linearly independent holomorphic solutions w_j , $j = 1, 2$, of the form

$$(3.5) \quad w_j(\xi) = \exp \{ (-1)^{j+1} \lambda_n \xi \} (1 + \varepsilon_j(\xi)) ,$$

and such that

$$(3.6) \quad |\varepsilon_j(\xi)| \leq \exp \left(\frac{1}{|\lambda_n|} \int_{s_j}^{\xi} |g_n(t)| |dt| \right) - 1, \quad \xi \in [a, b] \times [c, d] .$$

The integrals here are taken following the *progressive paths*, i.e. contours along which $\operatorname{Re}(\xi)$ is non-decreasing (for $j = 1$) or non-increasing (for $j = 2$). In the case of the rectangle we may take $s_1 = a + ic$ and $s_2 = b + id$, so that the whole rectangle is reachable by progressive paths. Taking advantage of the fact that the coefficients of the original equation (1.2) are polynomials, we can estimate the total variation

$$M(D) = \max_{\gamma} \int_{\gamma} |g_n(t)| |dt| ,$$

where γ is any horizontal or vertical segment in $[a, b] \times [c, d]$. This yields the following result:

Proposition 3.1. *With the assumptions above, let \tilde{D} be an Euclidean rectangle such that $\xi_n(D) \subset \tilde{D}$ and ξ_n can be continued holomorphically to \tilde{D} in such a way that $\xi_n(\mathcal{A} \cup \mathcal{V}_n) \cap \tilde{D} = \emptyset$. Let $d_n = \operatorname{dist}(\partial \tilde{D}, \xi_n(D))$ be the Euclidean distance from $\xi_n(D)$ to the boundary of \tilde{D} . Then we can replace the estimates (3.6) by*

$$(3.7) \quad |\varepsilon_j(\xi)| \leq \exp \left(\frac{K}{d_n |\lambda_n|} \right) - 1, \quad \xi \in \mathcal{R} .$$

In a ϖ_n -rectangle D we can select a single valued branch of the function H_n introduced in (2.5). Then a direct consequence of the proposition above is

Corollary 3.2. *Let D be a ϖ -rectangle. Then a general solution of (1.2) in D has the form*

$$(3.8) \quad y(z) = H_n(z) [\kappa_1 \zeta_n^{\lambda_n}(z) (1 + \varepsilon_1(z)) + \kappa_2 \zeta_n^{-\lambda_n}(z) (1 + \varepsilon_2(z))] ,$$

with ζ_n defined in (2.2). We have

$$(3.9) \quad |\varepsilon_j(z)| \leq \exp \left(\frac{K}{d_n |\lambda_n|} \right) - 1, \quad z \in D ,$$

where d_n is the Euclidean distance defined in Proposition 3.1.

3.4. Global asymptotic formula away from zeros. The result above shows that if we stay away from the singularities $\mathcal{A} \cup \mathcal{V}_n$, we can control the errors in the WKB approximation uniformly. This motivates the following definition. For $\varepsilon > 0$ and $t \in \mathbb{C}$ and subset $K \subset \mathbb{C}$ we denote

$$D_{\varepsilon}(t) \stackrel{\text{def}}{=} \{z \in \mathbb{C} : |z - t| < \varepsilon\}, \quad D_{\varepsilon}(K) \stackrel{\text{def}}{=} \bigcup_{t \in K} D(t, \varepsilon) ,$$

and

$$D_{n, \varepsilon} \stackrel{\text{def}}{=} D_{\varepsilon}(\mathcal{A} \cup \mathcal{V}_n) .$$

Let D be a ϖ_n -rectangle in $\mathbb{C} \setminus D_{n,\varepsilon}$ containing infinity. We can select the branch of ζ_n and the initial point z_0 in (2.2) in such a way that $|\zeta_n(z)| \geq 1 + \delta$ for a $\delta > 0$, and there is a constant κ such that

$$(3.10) \quad Q_n(z) = H_n(z) \left[\zeta_n^{\lambda_n}(z) (1 + \varepsilon_1(z)) + \kappa \zeta_n^{-\lambda_n}(z) (1 + \varepsilon_2(z)) \right].$$

It follows that if we continue all functions analytically in the maximal domain containing D such that $|\zeta_n(z)| \geq 1 + \delta$, then formula (3.10) still holds.

Next, standard arguments show that in a domain D , if $Q_n \neq 0$ for all sufficiently large n , then there exists a dominant term in (3.8):

Lemma 3.3. *Let $D = \{z \in \mathbb{C} : a < \operatorname{Re}(\xi_n) < b, c < \operatorname{Im}(\xi_n) < d\}$ be a simply connected ϖ_n -rectangle in $\mathbb{C} \setminus D_{n,\varepsilon}$, with*

$$d - c > \frac{\pi}{|\lambda_n|}.$$

If Ω does not contain zeros of a solution y of (1.2) then for any choice of the lower limit of integration and with an appropriate choice of the branch of the square root in (2.2) we have

$$(3.11) \quad y(z) = H_n(z) \zeta_n^{\lambda_n}(z) (1 + \varepsilon_n(z)).$$

Moreover, there exists a constant $M = M(\varepsilon, D)$, independent of n , such that for any ϖ_n -rectangle $D' \subset D$,

$$(3.12) \quad |\varepsilon_n(z)| \leq \frac{M}{|\lambda_n|}, \quad t \in D'.$$

Now let us observe what happens if we have a union of two ϖ_n -rectangles, $D^{(1)}$ and $D^{(2)}$, such that $D^{(1)} \cap D^{(2)}$ has an interior point, z_0 . Assume first that we take this z_0 as the lower limit of integration in (2.2) and (2.5). From Lemma 3.3 it follows that if $y \neq 0$,

$$y(z) = \kappa^{(j)} H_n(z) \zeta_n^{\lambda_n}(z) \left(1 + \varepsilon_n^{(j)}(z) \right), \quad z \in D^{(j)}, \quad j = 1, 2,$$

for certain non zero constants $\kappa^{(1)}$ and $\kappa^{(2)}$. Evaluating at $z = z_0$ we get

$$\kappa^{(1)} \left(1 + \varepsilon_n^{(1)}(z_0) \right) = \kappa^{(2)} \left(1 + \varepsilon_n^{(2)}(z_0) \right), \quad z \in D^{(1)} \cap D^{(2)},$$

so that,

$$\kappa^{(2)} = \kappa^{(1)} \frac{1 + \varepsilon_n^{(1)}(z_0)}{1 + \varepsilon_n^{(2)}(z_0)} = \kappa^{(1)} (1 + \varepsilon_n(z)) , \quad |\varepsilon_n(z)| \leq \frac{M'}{|\lambda_n|}.$$

In the intersection (neighborhood of z_0) both expressions for y should match, hence we have that in $D^{(1)} \cap D^{(2)}$,

$$\zeta_n^{\lambda_n}(z) \left(1 + \varepsilon_n^{(1)}(z) \right) = \zeta_n^{\pm \lambda_n}(z) (1 + \varepsilon_n(z)) \left(1 + \varepsilon_n^{(2)}(z) \right).$$

Since $D^{(1)} \cap D^{(2)}$ is an open set, necessarily the same branch of ζ_n has been taken in both sides of the previous identity. This argument shows that Lemma 3.3 is valid in the union $D^{(1)} \cup D^{(2)}$, eventually with a different constant in the right hand side of (3.12).

Now, modifying the lower limit of integration only changes the normalization constant κ in (3.11). That means that in the representation (3.11) we can choose

any lower bound $z_0 \in D^{(1)} \cup D^{(2)}$, as long as we understand the function ζ_n as the analytic continuation along a path in $D^{(1)} \cup D^{(2)}$.

The discussion above motivates the following definition:

Definition 3.4. Let $D^{(j)}$, $j = 1, \dots, k$, be a finite set of ϖ_n -rectangles with ϖ_n -height greater than $|\lambda_n|$ and bounded ϖ_n -perimeters. If each intersection $D^{(j)} \cap D^{(j+1)}$, $j = 1, \dots, k-1$, has an interior point, then their union $\cup_{j=1}^k D^{(j)}$ is called a finite ϖ_n -chain (see Figure 1).

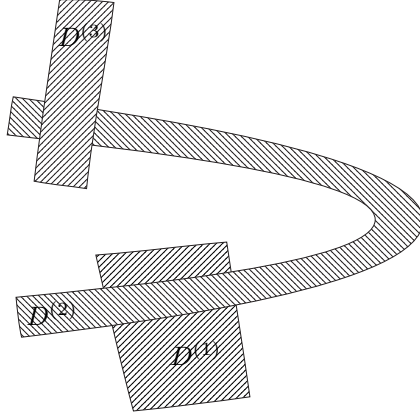


FIGURE 1. A ϖ_n -chain.

Thus, we have proved the following result:

Lemma 3.5. Let $\tilde{D} \subset \mathbb{C} \setminus D_{n,\varepsilon}$ be a domain, and $D = \cup_{j=1}^k D^{(j)} \subset \tilde{D}$ a finite and simply connected ϖ_n -chain in \tilde{D} . If a solution y of (1.2) does not vanish in D then for any choice of the lower limit of integration in D , and with an appropriate choice of the branch of the square root in (2.2) there exists a constant C such that (3.11)–(3.12) holds in D . The branch of ζ_n is obtained by analytic continuation in D .

In other words, asymptotic representation can be continued along the chains of ϖ_n -rectangles, as long as we stay away from the zeros of the solution and of the sets \mathcal{A} and \mathcal{V}_n .

3.5. Zeros of a solution of the differential equation. Assume now that y is a nontrivial solution of (1.2) and $z_0 \in \mathbb{C} \setminus D_{n,\varepsilon}$ is a zero of y . Let Ω be a maximal simply connected ϖ_n -rectangle in $\mathbb{C} \setminus D_{n,\varepsilon}$ containing z_0 , and let γ be the vertical trajectory in γ passing through z_0 . By assumption, z_0 is a regular point of the quadratic differential ϖ_n , so that γ is well defined.

In Ω the expression (3.8) is valid. In particular,

$$(3.13) \quad y(z) = 0 \quad \Leftrightarrow \quad 2\lambda_n \xi_n(z) = \log \left(-\frac{\kappa_2}{\kappa_1} \frac{1 + \varepsilon_1(z)}{1 + \varepsilon_2(z)} \right) \quad \text{mod } (2\pi i).$$

Hence, from the assumption $y(z_0) = 0$ it follows that

$$2\lambda_n \xi_n(z_0) = \log \left(-\frac{\kappa_2}{\kappa_1} \frac{1 + \varepsilon_1(z_0)}{1 + \varepsilon_2(z_0)} \right) \quad \text{mod } (2\pi i).$$

Assume that z_1 is another point on γ that satisfies the following condition:

$$(3.14) \quad \frac{\lambda_n}{\pi} \int_{z_0}^{z_1} |d\xi_n(t)| = \lambda_n \|\gamma(z_0, z_1)\|_{\varpi_n} \in \mathbb{N},$$

where $\gamma(z_0, z_1)$ is the arc of γ joining z_0 and z_1 .

Definition 3.6. Let $z_0 \in \mathbb{C} \setminus D_{n,\varepsilon}$, and let γ be the largest connected vertical arc in $\mathbb{C} \setminus D_{n,\varepsilon}$ of the quadratic differential ϖ_n passing thorough z_0 . The set

$$\omega = \bigcup_j \omega_j$$

is a *necklace* in $\mathbb{C} \setminus D_{n,\varepsilon}$ corresponding to z_0 if all ω_j , called *beads*, are ϖ_n -rectangles of the form

$$\omega_j = \{z \in \Omega : |\operatorname{Re}(\xi_n(z) - \xi_n(z_j))| < \delta, |\operatorname{Im}(\xi_n(z) - \xi_n(z_j))| < \delta\},$$

where each z_j satisfies condition (3.14), there exists a constant M such that $\lambda_n^2 \delta \leq M$, and all $\omega_j \subset \mathbb{C} \setminus D_{n,\varepsilon}$. The vertical arc γ is the *string* of the necklace.

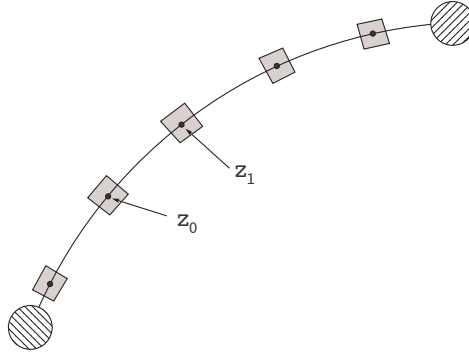


FIGURE 2. Necklace corresponding to z_0 .

A direct consequence of Rouché's theorem is the following statement:

Lemma 3.7. Let $z_0 \in \mathbb{C} \setminus D_{n,\varepsilon}$ be a zero of a nontrivial solution y of (1.2), and let ω be a necklace in $\mathbb{C} \setminus D_{n,\varepsilon}$ corresponding to z_0 . Then each bead $w_j \subset \mathbb{C} \setminus D_{n,\varepsilon}$ contains one and only one zero of y .

Furthermore, let Ω be a simply connected ϖ_n -rectangle in $\mathbb{C} \setminus D_{n,\varepsilon}$, and $z_0 \in \Omega$ be a zero of a nontrivial solution y of (1.2). If ω is the necklace in $\mathbb{C} \setminus D_{n,\varepsilon}$ corresponding to z_0 , then $y(z) \neq 0$ for $z \in \Omega \setminus \omega$.

These results are applicable to any solution of (1.2). Now we concentrate on the n -th degree Heine-Stieltjes polynomial. The main difference is that we know that it has exactly n zeros, with account of multiplicity. This and Lemma 3.7 immediately yield the following

Proposition 3.8. Assume that $z_0 \in \mathbb{C} \setminus D_{n,\varepsilon}$ is a zero of a Heine-Stieltjes polynomial Q_n , and let γ be the largest connected vertical arc in $\mathbb{C} \setminus D_{n,\varepsilon}$ of the quadratic differential ϖ_n passing thorough z_0 . Then

$$\|\gamma\|_{\varpi_n} \leq \frac{n+2}{\lambda_n}.$$

In particular, every string of a zero-carrying necklace has a finite ϖ_n -length, and starts and ends at $D_{n,\varepsilon}$, or is a closed curve.

Indeed, by Lemma 3.7, every necklace carries at least $\lambda_n \|\gamma\|_{\varphi_n} - 2$ zeros of Q_n , where γ is the string of the necklace.

Summarizing, if Q_n is a Heine-Stieltjes polynomial, there exist a finite number of zero-carrying necklaces with the corresponding strings $\gamma_{n,j}$ such that all the zeros belong to the union of these necklaces or lie in $D_{n,\varepsilon}$.

4. Asymptotics of Van Vleck polynomials

As it is clear from Theorem 2.1, the zeros of the Van Vleck polynomial V_n are the key parameters in the asymptotic expression for the Heine-Stieltjes polynomials. Hence, as a next step we derive a set of equations that will characterize their positions. We start with a formal argument and postpone to next section a more detailed discussion about consistency and meaning of these equations.

Theorem 4.1. *With the assumptions and notations above, let $y = Q_n$ be a Heine-Stieltjes polynomial (solution of (1.2)) corresponding to V_n , and let $\gamma_{k,n}$ be the set of arcs defined in Theorem 2.1. If there exists an $\varepsilon > 0$ such that all arcs $\gamma_{k,n}$ are disjoint and the ϖ_n -distance between them is $> \varepsilon$, then the following system of equations is satisfied:*

$$(4.1) \quad \frac{1}{\pi i} \int_{\gamma_{k,n}} \sqrt{\frac{V_n(t)}{A(t)}} dt = \frac{m_k}{\lambda_n} + \frac{\rho_k/2 + \delta_{k,n}}{\lambda_n}, \quad k = 1, \dots, p-1, \quad m_k \in \mathbb{M}_n,$$

where the index set \mathbb{M}_n is a finite subset of $\mathbb{N} \cup \{0\}$. If for $\gamma_{k,n}$ we denote by η_k the set of endpoints of $\gamma_{k,n}$ that belong to \mathcal{A} , then

$$\rho_k \stackrel{\text{def}}{=} \sum_{a \in \eta_k} \frac{B(a)}{A'(a)} + 1 - \text{card}(\eta_k).$$

Furthermore, there exists a constant $C = C(\varepsilon)$ such that

$$|\delta_{k,n}| \leq \frac{C}{n}.$$

This theorem is a simple consequence of the asymptotic formula (2.4), which is valid in a neighborhood of $\gamma_{k,n}$. Observing its increment along a closed Jordan curve $\tilde{\gamma}_{k,n}$ encircling $\gamma_{k,n}$ in the positive direction and using the argument principle we get the formulas above.

The system of equation (4.1) consists of $p-1$ equations (since we have fixed the residue at infinity, equation $k = p$ is dependent from the other $p-1$ ones). More exactly, (4.1) presents a collection of systems of equations. In order to define it completely we need to specify:

- (a) the combinatorics: points from $\mathcal{A} \cup \mathcal{V}_n$ are arranged in p pairs (endpoints of $\gamma_{k,n}$).
- (b) the homotopic types of curves $\gamma_{k,n}$ (once the combinatorics is fixed); and
- (c) the range of values the integer parameters m_k in the right hand side of (4.1) may take.

Observe that by (4.1),

$$\frac{1}{\pi i} \int_{\gamma_{k,n}} \sqrt{\frac{V_n(t)}{A(t)}} dt = \frac{m_k}{\lambda_n} + \frac{\rho_k/2}{\lambda_n} + \mathcal{O}\left(\frac{1}{n^2}\right), \quad k = 1, \dots, p-1, \quad m_k \in \mathbb{N}.$$

This system allows eventually to find the position of “almost all” zeros of Van Vleck polynomials, with an error smaller than the distance between zeros. In most part of the range this error is $\mathcal{O}(n^{-2})$. A complete analysis however is cumbersome and contains a combination of analytic, geometric and combinatorial arguments. We restrict our presentation here to the case $p = 2$ (three points), which at least has a trivial combinatorics and a rather simple geometry, with additional remarks on the case $p = 3$.

4.1. Case of $p = 2$. We have $\mathcal{A} = \{a_0, a_1, a_2\}$ (in general, non-collinear) and want to discuss the issues (a)-(c) raised above in order to define completely the set of equations determining the position of the zero of the Van Vleck polynomial. The first ingredient we need is the solution of the classical minimal capacity problem posed in the class of all continua in \mathbb{C} containing \mathcal{A} (the Chebotarev’s problem). It was proved by Grötzsch [3] and Lavrentiev [6, 7] that there exists a unique $\Gamma^* = \Gamma^*(\mathcal{A})$ satisfying

$$(4.2) \quad \text{cap}(\Gamma^*) = \min\{\text{cap}(F) : F \text{ a continuum containing } \mathcal{A}\},$$

where $\text{cap}(\cdot)$ denotes the logarithmic capacity. If a_j ’s are not collinear, then Γ^* is a union of three arcs, γ_0^* , γ_1^* , and γ_2^* , connecting a point $v^* = v^*(\mathcal{A})$ with points a_j ’s, respectively (see also [5] and [12]). We call this Γ^* the *Chebotarev’s compact* or *Chebotarev’s continuum* corresponding to \mathcal{A} , and point v^* is the *Chebotarev’s center* of the set \mathcal{A} . If we define

$$(4.3) \quad M_j = \frac{1}{\pi} \int_{\gamma_k} \left| \frac{t - v^*}{A(t)} \right|^{1/2} |dt|, \quad k = 0, 1, 2,$$

then

$$M_0 + M_1 + M_2 = 1.$$

Observe that each M_j is the ϖ^* -length of γ_j^* , where

$$\varpi^* \stackrel{\text{def}}{=} \frac{v^* - z}{A(z)} (dz)^2.$$

Next we define three analytic functions (elements) $w_k(v)$, $k = 0, 1, 2$. We describe first w_0 as a germ of an analytic function at $v = a_0$, which allows unlimited analytic continuation to $\mathbb{C} \setminus \mathcal{A}$. For all v in a sufficiently small neighborhood of a_0 let Δ_0 be the segment $[a_0, v]$, and $\Delta_1 = [a_1, a_2]$. Denote $\Omega \stackrel{\text{def}}{=} \mathbb{C} \setminus (\Delta_0 \cup \Delta_1)$. With $R(z) \stackrel{\text{def}}{=} (z - v^*)/A(z)$ we consider in Ω the single-valued branch of \sqrt{R} given by the asymptotic condition $\lim_{z \rightarrow \infty} z\sqrt{R(z)} = 1$, and define

$$w_0(v) = \frac{1}{\pi i} \int_{\Delta_0} \sqrt{R(t)} dt = \frac{1}{2\pi i} \oint_{\partial\Delta_0} \sqrt{R(t)} dt,$$

where $\partial\Delta_0$ is the doubly-connected component of $\partial\Omega$ (the Carathéodory boundary of Ω), with the boundary values of the branch of \sqrt{R} specified above. It is clear that w_0 is analytic at $v = a_0$, and may be continued as a multi-valued analytic function to $\mathbb{C} \setminus \mathcal{A}$. Similarly, we define w_1 and w_2 , starting from $v = a_1$ and $v = a_2$, respectively. The proof of the following statement can be found, e.g. in [8]:

Proposition 4.2. *For $k = 0, 1, 2$ there exist an analytic arc ℓ_k , connecting a_k with v^* , defined by*

$$\ell_k = \{v \in \mathbb{C} : w_k(v) \in [0, M_k]\}.$$

Furthermore, for $k = 0, 1, 2$, $w_k(v)$ is univalent in a neighborhood \mathcal{L}_k of ℓ_k (which is mapped by w_k onto a neighborhood of the interval $[0, M_k]$).

Now we are ready to formulate a theorem about the asymptotic location of the zeros of Van Vleck polynomials for $p = 2$.

Theorem 4.3. *Let $y = Q_n$ be a Heine-Stieltjes polynomial (solution of (1.2)) with $\mathcal{A} = \{a_0, a_1, a_2\}$, corresponding to $V_n(z) = z - v_n$. Assume that there exists an $\varepsilon > 0$ such that $\text{dist}(v_n, v^*) > \varepsilon$ (where v^* is the Chebotarev's center). Then for all sufficiently large n , there exist an index $k \in \{0, 1, 2\}$ and an integer value $0 \leq m_k \leq [nM_n]$ such that the following equation is satisfied:*

$$(4.4) \quad w_k(v_n) = \frac{m_k}{\lambda_n} + \frac{1}{\lambda_n} \frac{B(a_k)}{2A'(a_k)} + \frac{\delta_{k,n}}{\lambda_n},$$

and there exists a constant $C = C(\varepsilon)$ such that

$$|\delta_{k,n}| \leq \frac{C}{n}.$$

Obviously, this statement does not cover, roughly speaking, εn zeros out of $n + 1$ possible zeros of Van Vleck polynomials, but it provides an error estimate of order $C(\varepsilon)n^{-2}$. The exceptional set of “missed” Van Vleck zeros may be made smaller (up to a constant) for the price of relaxing the error estimate; also the technical details become more cumbersome.

4.2. Case of $p = 3$. Now we turn to the problems (a) and (b) related to the system (4.1), namely, we discuss the combinatorics and the homotopic type of curves $\gamma_{k,n}$ from Theorem 4.1; we leave the issue (c) of the range of values of m_k in the equations to the following Section.

Here we illustrate the situation considering the points a_j from \mathcal{A} at vertices of a rectangle whose height (vertical size) is smaller then length (horizontal size), see Fig. 3, with the results of some numerical experiments, as well as Fig. 4, where the corresponding Chebotarev's continuum is depicted.

In this case $V_n(z) = (z - v_{1,n})(z - v_{2,n})$; its zeros are determined completely by a system of two independent equations of the form (4.1).

We claim that in this situation there are exactly $8 + 1$ homotopically different groups of systems, see Figure 6. In eight of them we integrate along two arcs $\gamma_{k,n}$, each connecting a point from \mathcal{A} with a zero of V_n . In the remaining case (Figure 6, bottom right) one of the curves will connect both zeros of V_n .

For the case depicted in Figure 6, upper left, system (4.1) is written as

$$\begin{cases} w_1(c_1, c_2) = \frac{1}{\pi i} \int_{a_1}^{v_1} \sqrt{\frac{V_n(t)}{A(t)}} dt = \frac{m_{k,1}}{\lambda_n} + \frac{1}{\lambda_n} \frac{B(a_1)}{2A'(a_1)} + \frac{\delta_{k,n,1}}{\lambda_n}, \\ w_2(c_1, c_2) = \frac{1}{\pi i} \int_{a_2}^{v_2} \sqrt{\frac{V_n(t)}{A(t)}} dt = \frac{m_{j,2}}{\lambda_n} + \frac{1}{\lambda_n} \frac{B(a_2)}{2A'(a_2)} + \frac{\delta_{j,n,2}}{\lambda_n}, \end{cases}$$

where $(m_{k,1}, m_{j,2})$ is a pair of integers. We postpone the discussion about the possible range of variation of these constants until next section. In order to claim that

$$|\delta_{k,n,j}| \leq \frac{C}{n}$$

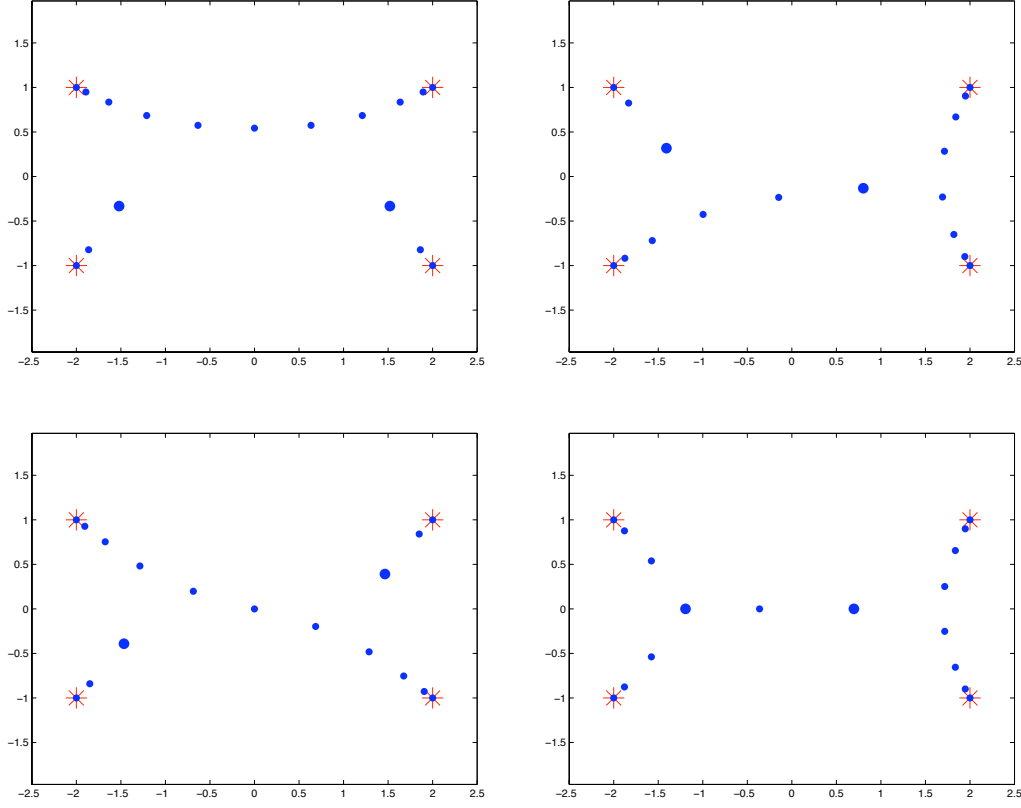


FIGURE 3. Zeros of Heine-Stieltjes (small dots) and of the corresponding Van Vleck polynomials (fat dots) for a_j 's at vertices of a rectangle ($p = 3$).

we need to leave aside again a fraction of zeros of Van Vleck polynomials (for which conditions of Theorem 2.1 are not satisfied), which is typically $o(n^2)$ (recall that the cardinality of the set of all Van Vleck polynomials corresponding to a Heine-Stieltjes polynomial of degree n is $\sigma(n) = \mathcal{O}(n^2)$).

5. The Van Vleck set (outline of the proof)

This section is essentially based on the results of our recent paper [8], and we will be regularly referring the reader to some parts of this work for further details.

We have seen that even in the simplest cases the geometry behind the set of equations (4.1) is quite involved; furthermore, we have not clarified yet the selection of the index set \mathbb{M}_n . However, the situation becomes much more clear when we take limit as $n \rightarrow \infty$.

We start by considering an arbitrary p for the price of omitting some non-essential details; a more thorough analysis will be carried out at the end for the particular (but non-trivial) case $p = 2$.

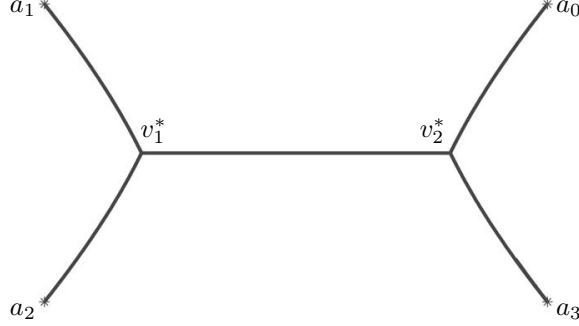


FIGURE 4. Chebotarev's compact corresponding to 4 points forming a rectangle.

5.1. Zero distribution of Heine-Stieltjes polynomials and critical measures. It is known (see e.g. [16]) that the zeros of the Heine-Stieltjes polynomials accumulate on the convex hull of \mathcal{A} . Hence, without loss of generality we may assume that there is a subsequence $\Lambda \subset \mathbb{N}$ such that the Van Vleck polynomials have a limit,

$$(5.1) \quad \lim_{n \in \Lambda} V_n(z) = V(z) = \prod_{j=1}^{p-1} (z - v_j).$$

This limit induces the following rational quadratic differential on the Riemann sphere,

$$(5.2) \quad \varpi = -\frac{V(z)}{A(z)} (dz)^2.$$

Let us consider the corresponding sequence of Heine-Stieltjes polynomials Q_n , $n \in \Lambda$; from Proposition 3.8 it follows that the zeros of Q_n 's lie asymptotically on critical trajectories of ϖ , and that these trajectories have the total ϖ -length 1. Further results can be obtained using the electrostatic interpretation of these zeros; it has been proved in [8] that any weak-* limit of the normalized zero-counting measure for Q_n 's is a continuous critical measure with respect to the set \mathcal{A} of fixed points on the plane. Let us recall the definition in its simplest form, sufficient for what follows.

With every (real-valued) Borel measure μ on \mathbb{C} we associate its (continuous) logarithmic energy

$$(5.3) \quad E(\mu) \stackrel{\text{def}}{=} \iint \log \frac{1}{|x - y|} d\mu(x) d\mu(y).$$

Any smooth complex-valued function h in the closure $\overline{\Omega}$ of a domain Ω containing \mathcal{A} generates a local variation of Ω by $z \mapsto z^t = z + t h(z)$, $t \in \mathbb{C}$, and consequently, a variation of sets $e \mapsto e^t \stackrel{\text{def}}{=} \{z^t : z \in e\}$, and (signed) measures: $\mu \mapsto \mu^t$, defined by $\mu^t(e^t) = \mu(e)$.

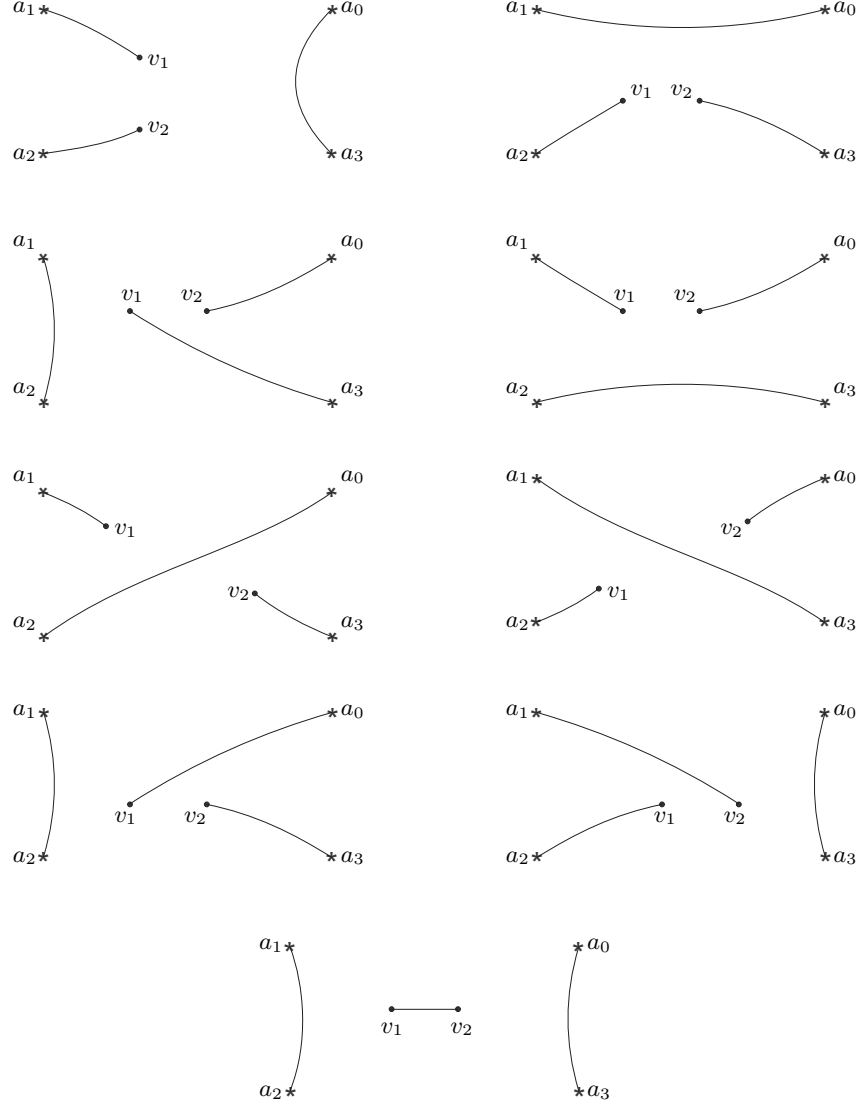


FIGURE 5. Possible homotopic classes of curves $\gamma_{n,k}$ for the case $p = 3$; compare with the examples depicted in Figure 3.

We say that a signed measure μ supported in Ω is a continuous \mathcal{A} -critical if for any h smooth in $\Omega \setminus \mathcal{A}$ such that $h|_{\mathcal{A}} \equiv 0$,

$$(5.4) \quad \frac{d}{dt} E(\mu^t) \Big|_{t=0} = \lim_{t \rightarrow 0} \frac{E(\mu^t) - E(\mu)}{t} = 0.$$

\mathcal{A} -critical measures appear quite frequently in many problems of approximation theory; any such a measure μ can be characterized in terms of its logarithmic

potential

$$U^\mu(z) \stackrel{\text{def}}{=} \int \log \frac{1}{|z-t|} d\mu(t)$$

as follows:

Lemma 5.1 ([8], Section 5.3). *The logarithmic potential of an \mathcal{A} -critical measure μ satisfies the following properties:*

- (i) *if $\text{supp}(\mu) = \gamma_1 \cup \dots \cup \gamma_s$, where γ_j are the connected components of $\text{supp}(\mu)$, then*

$$U^\mu(z) = w_j = \text{const}, \quad z \in \gamma_j, \quad j = 1, \dots, s.$$

- (ii) *at any regular point $z \in \text{supp}(\mu)$ (that is, such that locally at z , $\text{supp}(\mu)$ is a smooth Jordan arc),*

$$(5.5) \quad \frac{\partial U^\mu}{\partial n_+}(z) = \frac{\partial U^\mu}{\partial n_-}(z),$$

where n_\pm are the normal vectors to $\text{supp}(\mu)$ at z pointing in the opposite directions.

Additionally, if $z \in \text{supp}(\mu) \setminus \mathcal{A}$ is not regular, then

$$(5.6) \quad \text{grad } U^\mu(z) = 0.$$

Reciprocally, assume that a finite real measure μ , whose support $\text{supp}(\mu)$ consists of a union of a finite set of analytic arcs, $\text{supp}(\mu) = \gamma_1 \cup \dots \cup \gamma_s$, satisfies conditions (i) and (ii) above. Then μ is \mathcal{A} -critical.

In other words, (i) says that \mathcal{A} -critical measures are in fact equilibrium measures in a piece-wise constant external field, exhibiting additionally (see (ii)) the so-called S -property, introduced first by Stahl [17] and in a more general context, by Gonchar and Rakhmanov [2] (where it was used to establish the well-known “1/9”-conjecture in approximation theory). A first rigorous proof of the connection of a critical measure with the S -property appeared in [14].

With a polynomial $Q(z) = \prod_{k=1}^n (z - \zeta_k)$ of degree n we associate its (normalized) zero counting measure

$$\nu(Q) = \frac{1}{n} \sum_{k=1}^n \delta_{\zeta_k}.$$

The differential equation (1.2) is an expression of the fact that the zeros of Q_n sit in equilibrium (zeros of the gradient of the discrete total energy) in presence of an external field depending from the residues ρ_k in (1.4). The intensity of this field decays with n proportionally to $1/n$, so that in the limit we get

Proposition 5.2 ([8], Section 7). *Let $\nu_n = \nu(Q_n)$ be a zero-counting measure corresponding to a sequence of Heine-Stieltjes polynomials Q_n . Then any weak-* limit point μ of ν_n is a unit continuous \mathcal{A} -critical measure.*

In other words, weak-* limits of the normalized zero counting measures of Q_n ’s are unit positive \mathcal{A} -critical measures. The inverse inclusion (that any unit positive \mathcal{A} -critical measure is a weak-* limit of the normalized zero counting measures of Heine-Stieltjes polynomials) is also valid, but it cannot be established within the framework of this paper.

Continuous critical measures can be also related with rational quadratic differentials on $\overline{\mathbb{C}}$ by means of a variational argument. Our WKB analysis allows to establish this link, at least partially, more directly:

Theorem 5.3. *Let $y = Q_n$ be a Heine-Stieltjes polynomial (solution of (1.2)) corresponding to V_n , and assume that (5.1) holds. Then the quadratic differential*

$$\varpi = -\frac{V}{A}(z) dz^2$$

is closed (all its trajectories are either critical or closed), and there is a probability \mathcal{A} -critical measure μ on \mathbb{C} such that the normalized zero counting measures $\nu(Q_n)$ converge (in a weak- sense and along the subset Λ) to μ . The support $\Gamma = \text{supp}(\mu)$ consists of critical trajectories of ϖ , $\mathbb{C} \setminus \Gamma$ is connected, and we can fix the single valued branch of $\sqrt{V/A}$ there by $\lim_{z \rightarrow \infty} z \sqrt{V(z)/A(z)} = 1$. With this convention,*

$$(5.7) \quad \lim_n |Q_n(z)|^{1/n} = \exp \left(\text{Re} \int^z \sqrt{\frac{V}{A}}(t) dt \right)$$

locally uniformly in $\mathbb{C} \setminus \Gamma$, where a proper normalization of the integral in the right hand side is chosen, so that

$$\lim_{z \rightarrow \infty} \left(\text{Re} \int^z \sqrt{\frac{V}{A}}(t) dt - \log |z| \right) = 0.$$

The analysis above yields the following addendum, that we state using the notation and assumptions of Theorems 2.1 and 5.3:

Proposition 5.4. *The support $\Gamma = \text{supp}(\mu)$ is comprised of p analytic arcs γ_k , such that*

$$\lim_{n \in \Lambda} \Gamma_n = \Gamma.$$

5.2. Structure of the family of positive critical measures. Theorem 5.3 is essential for understanding the asymptotics of Heine-Stieltjes polynomials, although this result is in a certain sense implicit, since it depends on the limit V of the Van Vleck polynomials V_n , that constitute therefore the main parameters of the problem. We must complement this description with the study of the set of all possible limits V .

To this end it is convenient to consider a correspondence between closed quadratic differentials and \mathcal{A} -critical measures, independently of their origin. In general, this is not a one-to-one correspondence, since many critical measures may correspond to the same closed quadratic differential. It was pointed out in [8] that the bijection between closed quadratic differentials of the form (5.2) and signed \mathcal{A} -critical measures is restored if we restrict ourselves to signed measure with a connected complement of the support. Indeed, any such a measure is supported on a finite union of analytic arcs, and these arcs are necessarily critical trajectories of a quadratic differential like in (5.2). Reciprocally, from (5.2) we can construct explicitly (using the Sokhotsky-Plemelj formulas) a measure μ that satisfies (5.4).

Thus, any closed quadratic differential uniquely generates an \mathcal{A} -critical measure μ with a connected complement to its support. In general, such a μ is a signed measure, and for our purposes one has to select those quadratic differentials associated with positive measures. One of the ways to solve this problem goes through

describing the global structure of the trajectories of closed rational quadratic differentials with fixed denominators on the Riemann sphere, and the corresponding parameters (numerators), and later extracting from this set the subset giving rise to positive unit \mathcal{A} -critical measures.

This way is implemented in for $p = 2$ at the end of this paper. For an arbitrary p such an investigation would be more difficult to carry out, in particular due to the heavy combinatorics involved. Below we present a direct construction from [8], based on the v -local coordinates that we introduce in the space of the \mathcal{A} -critical measures, as well as on the topological structure of the Chebotarev's compact for \mathcal{A} . We start by introducing the v -coordinates.

Let us recall the notation. We have the fixed set $\mathcal{A} = \{a_0, a_1, \dots, a_p\}$ of distinct points on \mathbb{C} , $A(z) = \prod_{j=0}^p (z - a_j)$. For any signed \mathcal{A} -critical measure μ there exists a rational quadratic differential ϖ on the Riemann sphere $\overline{\mathbb{C}}$ given by

$$(5.8) \quad \varpi(z) = -R(z) (dz)^2, \quad R(z) \stackrel{\text{def}}{=} \frac{V(z)}{A(z)}, \quad V(z) \stackrel{\text{def}}{=} \prod_{j=1}^{p-1} (z - v_j),$$

such that $\text{supp}(\mu) = \Gamma_\mu = \Gamma = \gamma_1 \cup \dots \cup \gamma_p$ is a union of trajectories of ϖ .

We begin by introducing the local coordinates under the assumption that measures are in general position; this notion of genericity (as opposed to some more special or coincidental cases that are possible) means in our context that Γ is comprised of exactly p disjoint arcs. In consequence, zeros v_j 's of V are simple, and vector $\mathbf{v} \stackrel{\text{def}}{=} \{v_1, \dots, v_{p-1}\} \in \mathbb{C}^{p-1}$ can be used as a local coordinate. Define also

$$\mathcal{V} \stackrel{\text{def}}{=} \{\mathbf{v} : \varpi \text{ is closed}\}.$$

Next we use these local coordinates to introduce the period mapping for our quadratic differentials. The Carathéodory boundary of $\overline{\mathbb{C}} \setminus \Gamma$ consists of p components $\hat{\gamma}_k \stackrel{\text{def}}{=} \gamma_k^+ \cup \gamma_k^-$, with a positive orientation with respect to $\overline{\mathbb{C}} \setminus \Gamma$. We can consider $\hat{\gamma}_k$ as cycles in $\overline{\mathbb{C}} \setminus \Gamma$ enclosing the endpoints of γ_k . Part of \mathcal{R} over $\overline{\mathbb{C}} \setminus \Gamma$ splits into two disjoint sheets, so we may consider $\hat{\gamma}_k$ as cycles on \mathcal{R} .

Let us define

$$(5.9) \quad w_k(\mathbf{v}) = w_k(\mathbf{v}, \Gamma) \stackrel{\text{def}}{=} \frac{1}{2\pi i} \oint_{\hat{\gamma}_k} \sqrt{R(z)} dz, \quad k = 1, \dots, p,$$

where $\sqrt{R}|_{\hat{\gamma}_k}$ are the boundary values of the branch of \sqrt{R} in $\overline{\mathbb{C}} \setminus \Gamma$ defined by $\lim_{z \rightarrow \infty} z \sqrt{R(z)} = 1$. Clearly, the boundary values $(\sqrt{R})_{\pm}$ on γ_k^{\pm} are opposite in sign. Therefore, with any choice of orientation of γ_k and a proper choice of $\sqrt{R} = (\sqrt{R})_+$ on γ_k , we will have

$$(5.10) \quad w_k(\mathbf{v}) = \frac{1}{\pi i} \int_{\gamma_k} \sqrt{R(z)} dz, \quad k = 1, \dots, p.$$

By the Cauchy residue theorem we have that $w_1 + \dots + w_p = 1$ for any $\mathbf{v} \in \mathbb{C}^{p-1}$. Thus, we can restrict the mapping $\mathbf{v} \mapsto \mathbf{w}$ to $p - 1$ components of $\mathbf{w} \stackrel{\text{def}}{=} (w_1, \dots, w_{p-1}) \in \mathbb{C}^{p-1}$. In this way, we have defined the mapping

$$(5.11) \quad \mathcal{P}(\cdot, \Gamma) : \mathbb{C}^{p-1} \rightarrow \mathbb{C}^{p-1} \quad \text{such that} \quad \mathcal{P}(\mathbf{v}, \Gamma) = \mathbf{w}(\mathbf{v}, \Gamma).$$

Each component function $w_j(v_1, \dots, v_{p-1})$ is analytic in each coordinate v_k (even if v_k is at one of the endpoints of γ_k). Once defined by the integral in (5.10), this analytic germ allows an analytic continuation along any curve in $\mathbb{C} \setminus \mathcal{A}$. Arcs γ_k

are not an obstacle for the continuation since the integral in (5.10) depends only on the homotopic class of Γ in $\mathbb{C} \setminus (\mathcal{A} \cup \mathbf{v})$. The homotopy of Γ is a continuous modification of all components simultaneously in such a way that they remain disjoint in all intermediate positions. Under this assumption we can continuously modify the selected branch of \sqrt{R} in $\overline{\mathbb{C}} \setminus \Gamma$ along with the motion of Γ .

Proposition 5.5. *Mapping $w = \mathcal{P}(\mathbf{v})$ is locally invertible at any $\mathbf{v}(v_1, \dots, v_{p-1}) \in (\mathbb{C} \setminus \mathcal{A})^{p-1}$ with $v_i \neq v_j$ for $i \neq j$.*

Furthermore, we have

Proposition 5.6. *Let μ^0 be an \mathcal{A} -critical measure such that $\Gamma = \text{supp}(\mu^0)$ has connected components $\gamma_1^0, \dots, \gamma_p^0$, and $\overline{\mathbb{C}} \setminus \Gamma$ is connected. Let $\varpi_0 = R_0(z)(dz)^2$ be the quadratic differential associated with μ^0 , where $R_0 = V_0/A$, and $\mathbf{v}^0 = (v_1^0, \dots, v_{p-1}^0)$ is the vector of zeros of V_0 . Assume that ϖ_0 is in general position (that is, all v_k 's are pairwise distinct and disjoint with \mathcal{A}). Then for an $\varepsilon > 0$ and any $m_j \in \mathbb{R}$, $j \in \{1, \dots, p-1\}$, satisfying*

$$|m_j - \mu^0(\gamma_j^0)| < \varepsilon, \quad j = 1, \dots, p-1,$$

there exists a unique solution $\mathbf{v} \in \mathcal{V}$ of the system

$$(5.12) \quad w_j(\mathbf{v}, \Gamma) = m_j, \quad j = 1, \dots, p-1.$$

The quadratic differential $\varpi = -R(z)(dz)^2$, $R(z) = \prod_{k=1}^{p-1}(z - v_k)/A(z)$, is closed, and the associated \mathcal{A} -critical measure μ satisfies $\mu(\gamma_j) = m_j$, $j = 1, \dots, p-1$, where $\text{supp}(\mu) = \gamma_1 \cup \dots \cup \gamma_p$ and $\text{supp}(\mu)$ is homotopic to $\text{supp}(\mu^0)$.

This Proposition introduces a topology in the set of \mathcal{A} -critical measures in general position. We call a *cell* any connected component of this topological space. A measure in general position preserves sign along any connected component of its support. This sign is subsequently preserved when we homotopically modify the measure within its cell. In particular, if μ is a positive \mathcal{A} -critical measure, then all measures in the same cell with μ are positive.

In v -coordinates a cell $G = G(\Gamma)$ in \mathcal{V} is a subspace of $\mathbb{C}^{p-1} \simeq \mathbb{R}^{2p-2}$, which is a manifold of the real dimension $p-1$, defined by $p-1$ real equations of the form

$$(5.13) \quad \text{Im } w_j(v) = \text{Im} \left(\frac{1}{2\pi i} \int_{\hat{\gamma}_j} \sqrt{R(t)} dt \right) = 0, \quad j = 1, \dots, p-1,$$

where $\hat{\Gamma} = \hat{\gamma}_1 \cup \dots \cup \hat{\gamma}_{p-1} \cup \hat{\gamma}_p$ is a union of Jordan contours $\hat{\gamma}_k$ on \mathcal{R} (double arcs) depending on v , but mutually homotopically equivalent for values of v from the same cell. Practically, any $v_0 \in G$ has a neighborhood of v satisfying (5.13) with constant $\hat{\gamma}_k$'s.

Next, let \mathcal{V}_+ be the set of positive \mathcal{A} -critical measures; it will be comprised of 3^{p-1} cells. Below we introduce these cells by means of the Chebotarev's continuum $\Gamma^* = \Gamma^*(\mathcal{A})$, associated with $\mathcal{A} = \{a_0, \dots, a_p\}$ (that is, the continuum solving the extremal problem (4.2)). It constitutes the “universal boundary”, i.e., the intersection of all boundary points of all cells. In general, isolating a cell from its boundary point is a formidable task. We take advantage here that we are interested only in those that yield positive \mathcal{A} -critical measures, that are characterized by the trivial homotopy.

The Chebotarev's continuum consists of critical trajectories of $\varpi = -R^*dz^2$, $R^* = V^*/A$. Assume that the set \mathcal{A} satisfies the condition that A and V^* do not

have common zeros and $V^*(z) = \prod_{k=1}^{p-1} (z - v_k^*)$ does not have multiple zeros. Thus, the critical set $\mathcal{A}^* = \mathcal{A} \cup \{v_1^*, \dots, v_{p-1}^*\}$ consists of $2p$ different points, and Γ^* is comprised of $2p - 1$ arcs, that are critical trajectories of ϖ . Each trajectory joins two different points from \mathcal{A}^* .

Now we come to the procedure of selection of combinatorial (rather than homotopic) types of cells; once the combinatorial type is fixed, the homotopic one will be determined from the Chebotarev's continuum, as described next. Each zero $v^* = v_k^*$, $k = 1, \dots, p - 1$, is connected by component arcs of Γ^* with three other points, say $a_1^*, a_2^*, a_3^* \in \mathcal{A}^*$. We select one of these three arcs (for definiteness, $[v^*, a_1^*]$) and join two other arcs to make a single arc $[a_2^*, a_3^*]$, bypassing v^* (we think that the arc $[a_2^*, a_3^*]$ still follows the two arcs from Γ^* , but without touching v^* , instead passing infinitely close to it). This procedure, carried out at each zeros v_k^* of V^* , creates a compact set Γ , and consequently, a cell $G(\Gamma)$ of corresponding measures $\mu \in \widehat{\mathcal{V}}_+$.

The selection of Γ , and hence, of the cell $G(\Gamma)$, is made by choosing one of the three connections for each v_k^* ; there are 3^{p-1} ways to make the choice. Any choice splits Γ^* into p "disjoint" arcs $\Gamma^* = \gamma_1 \cup \dots \cup \gamma_p$; out of them we select $p - 1$ arcs (to make an homology basis for $\mathbb{C} \setminus \Gamma^*$) and then consider the corresponding cycles $\widehat{\gamma}_k$, as described above.

Finally, we describe the cell $G(\Gamma)$ in terms of the mapping \mathcal{P} . Let $w = w(\mathbf{v}) = \mathcal{P}(v, \widehat{\Gamma})$; the cell $G(\Gamma)$ is completely defined by the system

$$(5.14) \quad w_j(\mathbf{v}) = \mu_j \in \mathbb{R}_+, \quad j = 1, \dots, p - 1;$$

more precisely, there exists a domain $M(\Gamma) = \{(\mu_1, \dots, \mu_{p-1}) \in \mathbb{R}_+^{p-1}\}$ such that for any point $(\mu_1, \dots, \mu_{p-1}) \in M(\Gamma)$ system (5.14) has a unique solution $\mathbf{v} \in \mathbb{C}^{p-1}$. Moreover, the corresponding measure $\mu = \mu_v$ satisfies $\mu(\gamma_j) = \mu_j$, and $\text{supp}(\mu) = \gamma_1 \cup \dots \cup \gamma_p = \Gamma_v$ is homotopic to Γ .

Summarizing, a rough description of the set $\widehat{\mathcal{V}}_+$ of unit positive \mathcal{A} -critical measures may be made as follows. The set $\widehat{\mathcal{V}}_+$ is a union of 3^{p-1} of closed bounded cells $\overline{G}(\Gamma)$ ($\Gamma = \gamma_1 \cup \dots \cup \gamma_p$ may be selected in 3^{p-1} ways). The interior $G(\Gamma)$ of each cell consists of measures μ in general position with $\text{supp}(\mu)$ homotopic to Γ . Interiors of different cells are disjoint. Chebotarev's measure μ^* (Robin measure of Γ^*) is the only common point of all boundaries: $\mu^* = \bigcap_{\Gamma} \partial G(\Gamma)$. A graphical description of all these cells for $p = 3$, obtained from the construction just described, is contained in Figure 6.

5.3. Case $p = 2$. In order to clarify the construction above let us discuss in more detail the simplest (but far from trivial) case of $p = 2$.

Let us introduce the following set of the plane. For the quadratic differential

$$(5.15) \quad \varpi_v = \frac{v - z}{A(z)} dz^2$$

define

$$\mathcal{V} \stackrel{\text{def}}{=} \{v \in \mathbb{C} : \varpi_v \text{ is closed}\},$$

as well as the *Van Vleck set*

$$\mathcal{V}_+ \stackrel{\text{def}}{=} \{v \in \mathbb{C} : v \text{ is an accumulation point of the zeros of Van Vleck polynomials}\}.$$

A direct consequence of Theorem 5.3 is that

$$\mathcal{V}_+ \subset \mathcal{V}.$$

As it was mentioned for the general case, this inclusion is proper.

As in Subsection 4.1, we denote by v^* the Chebotarev's center of \mathcal{A} (the value of v in (5.15) such that ϖ_{v^*} corresponds to the Chebotarev's continuum for \mathcal{A}), and let $\mathcal{A}^* = \mathcal{A} \cup \{v^*\}$.

Theorem 5.7 ([8], Section 8.3). *The set \mathcal{V} is a union of a countable number of analytic arcs ℓ_k , $k \in \mathbb{Z}$, each connecting v^* and ∞ .*

Two arcs from \mathcal{V} are either identical or have v^ as the only finite common point. The homotopic type of the critical trajectories of ϖ_v in $\mathbb{C} \setminus \mathcal{A}$ remains invariant on each arc $\ell_k \setminus \mathcal{A}^*$.*

There are three distinguished arcs ℓ_k , $k \in \{0, 1, 2\}$, such that

- (i) ℓ_k connects v^* with infinity and passes through a_k ;
- (ii) for every $v \in \ell_k$ the homotopic class of trajectories of the closed quadratic differential ϖ_v is trivial.

In the terminology introduced in the previous subsection, this theorem says that each cell (connected component) in the topological space of \mathcal{A} -critical measures is homeomorphic to an analytic arc with endpoints either at v^* , \mathcal{A} or infinity. This settles the problem of identification of all closed quadratic differentials φ_v . Using the arguments described above we single out the differentials corresponding to *positive* measures:

Theorem 5.8. *Let ℓ_k , $k \in \{0, 1, 2\}$, be the distinguished arcs in \mathcal{V} described in Theorem 5.7. The set $\widehat{\mathcal{V}}_+$ is the union of the sub-arcs ℓ_k^+ of each ℓ_k , $k \in \{0, 1, 2\}$, connecting a_k with the Chebotarev's center v^* (and lying in the convex hull of \mathcal{A}).*

Furthermore, if $v \in \ell_k \cap \widehat{\mathcal{V}}_+$, $k \in \{0, 1, 2\}$, then there is a critical trajectory $\gamma(v)$ of μ_v connecting v with the pole a_k and such that

$$(5.16) \quad 0 \leq \mu_v(\gamma(v)) \leq M_k,$$

where M_k is defined in (4.3). In this case both critical trajectories of ϖ_v that constitute the support of μ_v are homotopic to a segment. The bijection $\mu_v(\gamma(v)) \leftrightarrow v$ is a parametrization of the set $\ell_k \cap \widehat{\mathcal{V}}_+$ by points of the interval $[0, m_k]$.

In other words, the only three cells corresponding to positive \mathcal{A} -critical measures are homeomorphic to three analytic arcs, joining the Chebotarev's center v^* with each pole $a_j \in \mathcal{A}$ (in fact, these arcs form a star homeotopic and quite close to the Chebotarev's continuum, although in general not coincident with it). If v travels such an arc, the boundary of the cell is reached when either endpoint of this arc is met. If we continue further along the same arc, a new cell is entered, corresponding to sign-changing \mathcal{A} -critical measures (see an illustration of the correspondence between the position of v on \mathcal{V} and the trajectories of ϖ_v in Figure 6).

Furthermore, by Theorem 5.8, the limit zero v of $V(z) = z - v$ satisfies an equation of the form

$$\frac{1}{\pi i} \int_{a_k}^v \sqrt{\frac{V(t)}{A(t)}} dt = \beta \in [0, M_k],$$

which is consistent with our construction in Subsection 4.1. We can go back now to the set of equations (4.1) and add the following result:

Theorem 5.9. *For any $\varepsilon > 0$, equations*

$$(5.17) \quad w_k(\tilde{v}_{j,k}) = \frac{j}{\lambda_n} + \frac{\rho_k/2}{\lambda_n}, \quad j = 0, \dots, M_k = [m_k n(1 - \varepsilon)],$$

with $k = 0, 1, 2$, uniquely define $\tilde{n} = (1 - \varepsilon)n \pm 1$ points $\tilde{v}_{j,k}$. There exists a constant $C = C(\varepsilon) > 0$ such that any Van Vleck zero lies in a C/n^2 distance from a point $\tilde{v}_{j,k}$.

In this way, each zero of V_n “belongs” to a $\mathcal{O}(n^{-2})$ neighborhood of one of the points defined by the equations above. We need to show additionally that (again, up to a small neighborhood of the Chebotarev’s center v^*) all neighborhoods of this form contain at least one Van Vleck zero (establishing in this form a bijection between points $\tilde{v}_{j,k}$ and the set \mathcal{V}_n). However, the techniques developed in this paper do not allow to complete this proof, so we formulate it as an open question:

Conjecture 5.10. *There exists a constant $C = C(\varepsilon) > 0$ such that to each point $\tilde{v}_{j,k}$ defined by equations (5.17) it corresponds at least one zero v of a Van Vleck polynomial such that $|v - \tilde{v}_{j,k}| \leq C/n^2$.*

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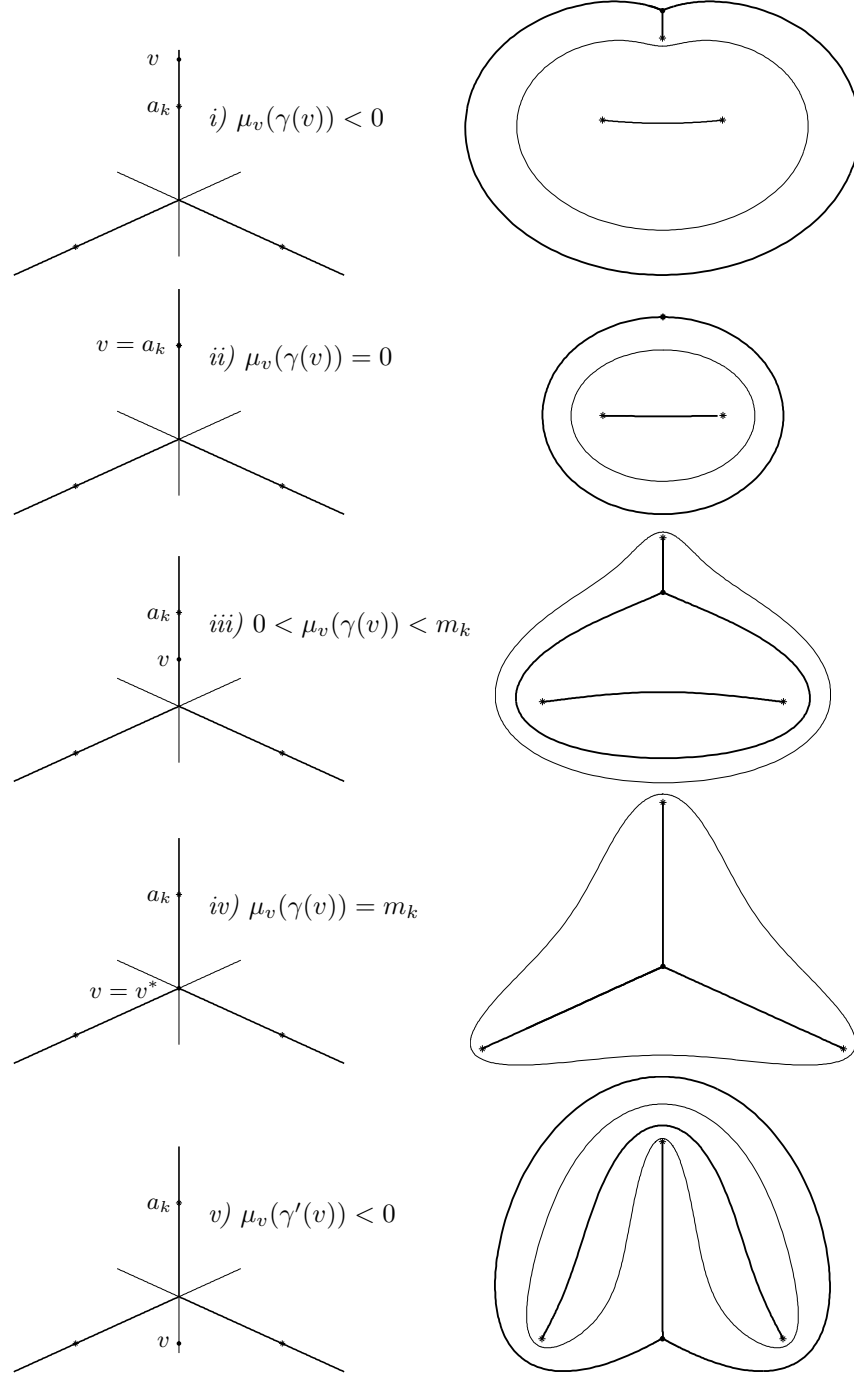


FIGURE 6. Position of v on $\ell_0 \cup \ell_1 \cup \ell_2$ (left) and the corresponding trajectories of the differential ϖ_v in (5.15).